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An Overdetermined Linear System

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1. INTRODUCTION

Consider a linear system of algebraic equations $Ax = b$, where $A = (a_{ij})$ is an $(n + 1) \times n$ real matrix, $b = (b_i)$ an $(n + 1)$ -component real vector and $x = (x_i)$ an n -component real vector. Clearly, the system is overdetermined if b is not in the range of A and it is this problem to which we seek a “best” approximate solution. The technique employed is to impose an abstract norm on R^{n+1} and then find $x \in R^n$ such that the error vector $\eta(x) = b - Ax$ is minimized with respect to this norm. We refer to this question of finding such a solution with respect to a given abstract norm as “Problem (P).” When we use the standard l^p -norm, we let $\xi_p = (\xi_i^p)$ denote an l^p -solution to $Ax = b$, i.e.,

$$\|\eta(\xi_p)\|_p = \min_{x \in R^n} \|\eta(x)\|_p.$$

(The least-squares solution and the Tchebychev solution correspond to the values $p = 2$ and $p = \infty$, respectively.)

A set of vectors in R^m is said to satisfy the Haar condition if every set of m of them is linearly independent. Let \mathcal{H} denote the set of all $(n + 1) \times n$ matrices whose rows satisfy the Haar condition. In this paper, we concern ourselves with the problem of finding the l^p -solution to $Ax = b$ when $A \in \mathcal{H}$. In Section 2, we give an explicit formula for the l^p -solution and also discuss how it may be expressed as a convex combination of the solutions to the $n \times n$ subsystems. In Section 3, we allow b to be a random vector and demonstrate how one may favor a particular norm when observing a minimum variance criterion. Two examples of such stochastic overdetermined systems are given. Finally, some preliminary results pertaining to the matrix A are presented in the Appendix.

2. EXPLICIT FORMS OF THE l^p -SOLUTION

Let A be an $(n+1) \times n$ matrix of rank n and A^T the transpose of A . Then $Ax = b$ has the unique l^2 -solution

$$\xi_2 = (A^T A)^{-1} A^T b.$$

Moreover, the l^2 -error vector is $s = \eta(\xi_2) = (I - A(A^T A)^{-1} A^T) b$ (see [1]).

The dual norm $\|\cdot\|^\sim$ of a norm $\|\cdot\|$ on R^{n+1} is defined by

$$\|u\|^\sim = \max_{\|v\|=1} (u, v),$$

where

$$(u, v) = \sum_{i=1}^{n+1} u_i v_i$$

is the standard Euclidean inner product on R^{n+1} . Furthermore, \tilde{u} is called a dual vector to $u \in R^{n+1}$ if $(\tilde{u}, u) = \|u\|^\sim$ and $\|\tilde{u}\| = 1$. From [4] we have the following

THEOREM (Sreedharan). *Let A be an $(n+1) \times n$ matrix of rank n and $s = \eta(\xi_2) = b - A\xi_2$ the l^2 -error vector. If $s = 0$, then ξ_2 is a solution of Problem (P). If $s \neq 0$, then*

$$Ax = b - ((b, s)/\|s\|^\sim) \tilde{s} \quad (2.1)$$

has a solution, and any solution of (2.1) is a solution of Problem (P).

From this point we shall only consider the case where R^{n+1} is equipped with the l^p -norm. Thus, the l^q -norm is the dual norm if $(1/p) + (1/q) = 1$. For $1 < p \leq \infty$, the dual vector of a nonzero vector $s = (s_i)$ is $\tilde{s}_p = (\tilde{s}_i^p)$, where

$$\tilde{s}_i^p = (|s_i|/\|s\|_q)^{q-1} \operatorname{sgn} s_i, \quad i = 1, 2, \dots, n+1. \quad (2.2)$$

For $p = 1$, we define $\rho = \{r_1, r_2, \dots, r_t\} = \{r: 1 \leq r \leq n+1 \text{ and } |s_r| = \max_i |s_i| = \|s\|_\infty\}$ and $\tilde{s}_1(r) = (\tilde{s}_i^1(r))$, where

$$\begin{aligned} \tilde{s}_i^1(r) &= \operatorname{sgn} s_r & \text{if } i = r, \\ &= 0 & \text{if } i \neq r. \end{aligned}$$

Then any convex combination of $\tilde{s}_1(r)$ with $r \in \rho$ is a dual vector of s for $p = 1$.

Clearly

$$\lim_{p \rightarrow 1^+} \tilde{s}_p = (1/t) \sum_{j=1}^t \tilde{s}_1(r_j) \quad (2.3)$$

is also a dual vector of s for $p = 1$.

With the given l^p -norm and $s \neq 0$, we may now write (2.1) as $Ax = b^{(p)}$ where

$$b^{(p)} = (b_i^{(p)}) = b - ((b, s)/\|s\|_q) \tilde{s}_p \quad (2.4)$$

and \tilde{s}_p is given by (2.2) for $1 < p \leq \infty$, while any convex combination of $\tilde{s}_1(r)$ with $r \in \rho$ can define \tilde{s}_1 . It is shown in [1, 2] that $Ax = b^{(p)}$ has a unique solution for each p , $1 \leq p \leq \infty$. Furthermore, each is a respective l^p -solution to $Ax = b$. The l^p -solution to $Ax = b$ is unique when $1 < p < \infty$ since the l^p -norm is strictly convex. From [2] we have that the l^1 -solution is unique if and only if ρ is a singleton set ($t = 1$), and the l^∞ -solution is unique if and only if the rows of A satisfy the Haar condition. If $Ax = b^{(p)}$ has a unique solution, then it can be found by solving $A^T Ax = A^T b^{(p)}$. Hence the l^p -solution can be obtained from

$$\xi_p = (A^T A)^{-1} A^T b^{(p)}. \quad (2.5)$$

Let A^k be the $n \times n$ matrix obtained from A by deleting the k th row, $D_k \equiv \det(A^k)$, the determinant of A^k ,

$$A_q \equiv \sum_{k=1}^{n+1} |D_k|^q$$

and

$$\sigma \equiv \sum_{k=1}^{n+1} (-1)^{k-1} b_k D_k.$$

Now $A \in \mathcal{H}$ implies that there exists a nonsingular $n \times n$ matrix P with $|\det(P)| = 1$ such that $AP = G$ where

$$G = \begin{pmatrix} g_1 & 0 & \cdots & 0 \\ 0 & g_2 & \cdots & 0 \\ \vdots & \vdots & & \vdots \\ 0 & 0 & \cdots & g_n \\ h_1 & h_2 & \cdots & h_n \end{pmatrix}$$

with $g_i \neq 0$ and $h_i \neq 0$, $i = 1, 2, \dots, n$. Let G^k be the matrix G with the k th row deleted. Then for $k = 1, 2, \dots, n+1$

$$A^k P = G^k$$

and

$$\hat{D}_k \equiv \det(G^k) = D_k \det(P).$$

Consider the $(n+1) \times (n+1)$ matrix formed by extending A in the following way:

$$(\alpha_{ij}) \equiv (A \mid \delta),$$

where $\delta \equiv ((-1)^{n+1-i} D_i)$ is an $(n+1)$ -component column vector. Define $m_{ij} \equiv (\text{cofactor of } \alpha_{ji})/\Delta_2$ for $i = 1, \dots, n, j = 1, \dots, n+1$, or more explicitly,

$$m_{ij} = \frac{(-1)^{i+j}}{\Delta_2} \begin{vmatrix} a_{11} & \cdots & a_{1,i-1} & a_{1,i+1} & \cdots & a_{1n} & (-1)^n D_1 \\ \vdots & & \vdots & \vdots & & \vdots & \vdots \\ a_{j-1,1} & \cdots & a_{j-1,i-1} & a_{j-1,i+1} & \cdots & a_{j-1,n} & (-1)^{n-j+2} D_{j-1} \\ a_{j+1,1} & \cdots & a_{j+1,i-1} & a_{j+1,i+1} & \cdots & a_{j+1,n} & (-1)^{n-j} D_{j+1} \\ \vdots & & \vdots & \vdots & & \vdots & \vdots \\ a_{n+1,1} & \cdots & a_{n+1,i-1} & a_{n+1,i+1} & \cdots & a_{n+1,n} & (-1)^0 D_{n+1} \end{vmatrix}. \quad (2.6)$$

THEOREM 2.1. Let $A \in \mathcal{H}$.

(i) If $\sigma = 0$, then $Ax = b$ is solvable and has the unique solution $x = Mb$ where $M = (m_{ij})$ is defined by (2.6).

(ii) If $\sigma \neq 0$, then an l^p -solution to $Ax = b$ is $\xi_p = B_p b$, where $B_p = (\beta_{ij}(p))$ is defined by

$$\begin{aligned} \beta_{ij}(p) &= m_{ij} + \frac{(-1)^j D_j}{\Delta_q} \sum_{k=1}^{n+1} (-1)^{k-1} \frac{m_{ik} |D_k|^q}{D_k}, \quad 1 < p \leq \infty, \\ &= m_{ij} + \frac{(-1)^j D_j}{\tau} \sum_{k \in \lambda} (-1)^{k-1} \frac{m_{ik}}{D_k}, \quad p = 1, \end{aligned}$$

$(1/p) + (1/q) = 1$, $\lambda = \{l_1, l_2, \dots, l_\tau\} = \{l: 1 \leq l \leq n+1 \text{ and } |D_l| = \max_i |D_i|\}$, and m_{ij} is given by (2.6). The l^1 -solution given here has the property that $\lim_{p \rightarrow 1^+} \xi_p = \xi_1$, hence ξ_p is continuous with respect to p for $1 \leq p \leq \infty$. Furthermore, ξ_p is the unique l^p -solution for $1 < p \leq \infty$ while the l^1 -solution is unique if and only if λ is a singleton set ($\tau = 1$).

Proof. Since $s = (I - A(A^T A)^{-1} A^T) b$, Theorem A.1 (cf. Appendix) implies that

$$s_i = (\sigma/\Delta_2)(-1)^{i-1} D_i, \quad i = 1, 2, \dots, n+1. \quad (2.7)$$

(i) $A \in \mathcal{H}$ implies that $\text{rank}(A) = n$ and thus A is one-to-one. Now b is in the range of A if and only if $s = 0$ which, in turn, is equivalent to $\sigma = 0$.

Hence $Ax = b$ is uniquely solvable in this case and it has the solution $x = (A^T A)^{-1} A^T b = Mb$.

(ii) For $\sigma \neq 0$, we see from (2.7) that s_i is proportional to D_i , so ρ and λ represent the same set. Furthermore,

$$(b, s) = \sigma^2 / \Delta_2. \quad (2.8)$$

From (2.2), (2.3), (2.7), (2.8), and (2.4) we obtain

$$b_j^{(p)} = b_j + \frac{\sigma}{\Delta_q} \frac{(-1)^j |D_j|^q}{D_j}, \quad 1 < p \leq \infty; \quad (2.9a)$$

$$\begin{aligned} b_j^{(1)} &= b_j + (\sigma(-1)^j / \tau D_j), & j \in \lambda, \\ &= b_j, & j \notin \lambda. \end{aligned} \quad (2.9b)$$

Theorem A.2 (cf. Appendix) and (2.5) yield

$$\xi_i^p = \sum_{j=1}^{n+1} m_{ij} b_j^{(p)}, \quad i = 1, 2, \dots, n; \quad 1 \leq p \leq \infty;$$

hence

$$\xi_i^p = \sum_{j=1}^{n+1} \beta_{ij}(p) b_j, \quad i = 1, 2, \dots, n, \quad 1 \leq p \leq \infty,$$

via (2.9) and the definition of σ . We see that $b_j^{(p)}$, $\beta_{ij}(p)$, and thus ξ_p are continuous with respect to p for $1 \leq p \leq \infty$ when \hat{s}_1 is defined by (2.3).

COROLLARY 2.1. *Let $A \in \mathcal{H}$, $\sigma \neq 0$, and $|D_i| = c$, $i = 1, 2, \dots, n+1$, for some c . Then ξ_p is independent of p for $1 \leq p \leq \infty$.*

Proof. This follows immediately from the definition of $\beta_{ij}(p)$.

Let us now consider the $n \times n$ subsystem of equations

$$A^k Z^k = b^k \quad (2.10)$$

for $k = 1, 2, \dots, n+1$, where b^k is the vector b with the k th element deleted. Then (2.10) has a unique solution Z^k since A^k is nonsingular ($D_k \neq 0$) by the Haar condition. We would like to establish a relation between the l^p -solution ξ_p and the Z^k 's. For each $k = 1, 2, \dots, n+1$, the $n \times n$ system

$$G^k W^k = b^k \quad (2.11)$$

also has a unique solution $W^k = (w_i^k)$ since $|\hat{D}_k| = |D_k| \neq 0$.

We note that

$$G^k = \begin{pmatrix} g_1 & \cdots & 0 & 0 & 0 & \cdots & 0 \\ \vdots & & \vdots & \vdots & \vdots & & \vdots \\ 0 & \cdots & g_{k-1} & 0 & 0 & \cdots & 0 \\ 0 & \cdots & 0 & 0 & g_{k+1} & \cdots & 0 \\ \vdots & & \vdots & \vdots & \vdots & & \vdots \\ 0 & & 0 & 0 & 0 & \cdots & g_n \\ h_1 & \cdots & h_{k-1} & h_k & h_{k+1} & \cdots & h_n \end{pmatrix}.$$

Hence the first $n - 1$ equations of the system (2.11) yield

$$w_i^k = b_i/g_i \quad \text{for } i \neq k. \quad (2.12a)$$

By Cramer's rule and through expanding the determinant in the numerator about the k th column, we obtain

$$w_k^k = \frac{1}{D_k} \sum_{j \neq k} (-1)^{k+j} b_j \left(-\frac{\hat{D}_j}{g_k} \right) = \frac{1}{D_k g_k} \sum_{j \neq k} (-1)^{k+j-1} b_j D_j. \quad (2.12b)$$

LEMMA 2.1. (i) If $\sigma = 0$ then $(G^T G)^{-1} G^T b = W \equiv (b_i/g_i)$.

(ii) If $\sigma \neq 0$ then

$$\begin{aligned} (G^T G)^{-1} G^T b^{(p)} &= (1/\Delta_a) \sum_{k=1}^{n+1} |D_k|^q W^k, & 1 < p \leq \infty, \\ &= (1/\tau) \sum_{k \in \lambda} W^k, & p = 1, \end{aligned}$$

where $b^{(p)}$, λ and τ were defined in Theorem 2.1.

Proof. Let $(G^T G)^{-1} G^T \equiv (\gamma_{ij})$. From the proof of Theorem A.1 we have, for $i = 1, 2, \dots, n, j = 1, 2, \dots, n+1$,

$$\begin{aligned} \gamma_{ij} &= \left(\sum_{k \neq i} D_k^2 \right) / (\Delta_2 g_i), & i = j, \\ &= (-1)^{i+j-1} D_i D_j / (\Delta_2 g_i), & i \neq j. \end{aligned} \quad (2.13)$$

(i) It follows from (2.13) and Lemma A.1 that $\sigma = 0$ implies

$$\sum_{j=1}^{n+1} \gamma_{ij} b_j = b_i/g_i, \quad i = 1, 2, \dots, n.$$

(ii) From (2.9a), we have

$$\sum_{j=1}^{n+1} \gamma_{ij} b_j^{(p)} = (1/\Delta_a) \sum_{k=1}^{n+1} |D_k|^q c_{ik}, \quad i = 1, 2, \dots, n, \quad 1 < p \leq \infty,$$

where

$$c_{ik} = \sum_{j=1}^{n+1} \gamma_{ij} b_j + (\sigma(-1)^k \gamma_{ik}/D_k).$$

From (2.13), the above equation yields

$$\begin{aligned} c_{ik} &= \frac{1}{D_k g_k} \sum_{j \neq k} (-1)^{k+j-1} b_j D_j, & i = k, \\ &= b_i/g_i, & i \neq k. \end{aligned}$$

Hence $c_{ik} = w_i^k$ (cf. (2.12)). This completes the proof for the case $1 < p \leq \infty$, while the result for $p = 1$ follows by letting $p \rightarrow 1^+$.

THEOREM 2.2. *Let $A \in \mathcal{H}$ and, for each $k = 1, 2, \dots, n+1$, let Z^k be the unique solution of $A^k Z^k = b^k$.*

(i) *If $\sigma = 0$, then $Ax = b$ is solvable and it has the unique solution $x = Z^k$, $k = 1, 2, \dots, n+1$ (all the Z^k 's being equal).*

(ii) *If $\sigma \neq 0$, then*

$$\xi_p = (1/\Delta_a) \sum_{k=1}^{n+1} |D_k|^q Z^k \quad (2.14)$$

is the unique l^p -solution to $Ax = b$ for $1 < p \leq \infty$. If λ and τ are defined as in Theorem 2.1, then

$$\xi_1 = (1/\tau) \sum_{k \in \lambda} Z^k \quad (2.15)$$

is an l^1 -solution to $Ax = b$ which is unique if and only if λ is a singleton set ($\tau = 1$). Furthermore, the l^p -solution provided here is continuous with respect to p for $1 \leq p \leq \infty$.

Proof. Since A^k is nonsingular, it follows from $A^k P = G^k$, (2.10), and (2.11) that

$$Z^k = PW^k, \quad k = 1, 2, \dots, n+1. \quad (2.16)$$

(i) From (2.12), $\sigma = 0$ implies $w_i^k = b_i/g_i$, $i = 1, 2, \dots, n$, which is independent of k . Hence $W^k = W \equiv (b_i/g_i)$, $k = 1, 2, \dots, n+1$, and $Z^k = PW^k = PW$, $k = 1, 2, \dots, n+1$. From Theorem 2.1, $Ax = b$ is uniquely solvable with solution $x = (A^T A)^{-1} A^T b$. Therefore (A.1) and Lemma 2.1 show that $x = PW = Z^k$.

(ii) From (2.5) and (A.1) we have

$$\xi_p = P(G^T G)^{-1} G^T b^{(p)}. \quad (2.17)$$

Hence (2.14) and (2.15) follow from (2.17), Lemma 2.1, and (2.16).

Remark. Theorem 2.2 demonstrates how an l^p -solution of an overdetermined system with $A \in \mathcal{H}$ may be expressed as a convex combination of the solutions of the $n \times n$ subsystems.

3. A STOCHASTIC OVERDETERMINED SYSTEM

Consider an overdetermined system $Ax = b$, where $A \in \mathcal{H}$ is a constant matrix but now $b = (b_i)$ is a random vector. Let $E(\cdot)$, $V(\cdot)$, and $\text{Cov}(\cdot, \cdot)$ denote the expected value, variance, and covariance operators, respectively. We assume $V(b_i) < \infty$ for $i = 1, 2, \dots, n+1$. From Theorem 2.1 we have

$$\xi_i^p = \sum_{j=1}^{n+1} \beta_{ij}(p) b_j, \quad i = 1, \dots, n, \quad 1 \leq p \leq \infty, \quad (3.1)$$

and hence

$$E(\xi_i^p) = \sum_{j=2}^{n+1} \beta_{ij}(p) E(b_j), \quad i = 1, \dots, n, \quad 1 \leq p \leq \infty. \quad (3.2)$$

We note that $E(\xi_p)$ is the l^p -solution to $Ax = E(b)$. We call $E(\xi_{\hat{p}})$ our "preferred approximate solution" to the stochastic overdetermined system if \hat{p} is selected according to the minimum variance condition

$$v(\hat{p}) = \min_{1 \leq p \leq \infty} v(p),$$

where

$$v(p) = \sum_{i=1}^n V(\xi_i^p)$$

is the trace of the covariance matrix. Note that $v(p)$ is invariant under orthogonal transformations and is equal to the sum of the eigenvalues. Then, in a sense, it measures the total spread of the random variables. If $|D_i| = c$ for $i = 1, 2, \dots, n+1$ then by Corollary 2.1 $v(p)$ is independent of p and thus any value of p , $1 \leq p \leq \infty$ will suffice for \hat{p} . Applying the variance operator to (3.1), we obtain

$$V(\xi_i^p) = \sum_{j=1}^{n+1} \beta_{ij}^2(p) V(b_j) + 2 \sum_{j < k} \beta_{ij}(p) \beta_{ik}(p) \text{Cov}(b_j, b_k).$$

Therefore

$$\nu(p) = \sum_{j=1}^{n+1} V(b_j) \sum_{i=1}^n \beta_{ij}^2(p) + 2 \sum_{j < k} \text{Cov}(b_j, b_k) \sum_{i=1}^n \beta_{ij}(p) \beta_{ik}(p). \quad (3.3)$$

We note that $\nu(p)$ is continuous for $1 \leq p \leq \infty$ and differentiable for $1 < p < \infty$.

Without loss of generality, we then assume that $|a_{11}| > |a_{21}|$ for the case $n = 1$. Here

$$M = (m_{11}, m_{12}) = \left(\frac{a_{11}}{a_{11}^2 + a_{21}^2}, \frac{a_{21}}{a_{11}^2 + a_{21}^2} \right),$$

$$\beta_{1j}(p) = \frac{|a_{j1}|^q}{a_{j1}(|a_{11}|^q + |a_{21}|^q)}, \quad 1 < p \leq \infty,$$

$$= (1 + (-1)^{j+1})/2a_{11}, \quad p = 1, \quad j = 1, 2,$$

and from (3.2) and (3.3)

$$E(\xi_p) = \frac{|a_{11}|^q \frac{E(b_1)}{a_{11}} + |a_{21}|^q \frac{E(b_2)}{a_{21}}}{|a_{11}|^q + |a_{21}|^q}, \quad 1 < p \leq \infty,$$

$$= E(b_1)/a_{11}, \quad p = 1,$$

$$\nu(p) = \frac{K_1}{(|a_{11}|^q + |a_{21}|^q)^2}, \quad 1 < p \leq \infty,$$

$$= V(b_1)/a_{11}^2, \quad p = 1,$$

where

$$K_1 = |a_{11}|^{2q-2} V(b_1) + |a_{21}|^{2q-2} V(b_2) + 2 \frac{|a_{11}a_{21}|^q}{a_{11}a_{21}} \text{Cov}(b_1, b_2).$$

If $V(b_1) = 0$ and $V(b_2) \neq 0$ we choose $\hat{p} = 1$, while if $V(b_1) \neq 0$ and $V(b_2) = 0$ we choose $\hat{p} = \infty$. For $V(b_1) \neq 0$ and $V(b_2) \neq 0$ we find $\lim_{p \rightarrow 1+} (d\nu/dp) = \lim_{p \rightarrow \infty} (d\nu/dp) = 0$. Furthermore, $d\nu/dp = 0$ at $p = p_0$, where

$$p_0 = 1, \quad \theta = 0,$$

$$= \infty, \quad \theta = 1, \quad (3.4)$$

$$= 1 + \frac{\ln(|a_{11}|/|a_{21}|)}{\ln|\theta|}, \quad \text{otherwise,}$$

with

$$\theta = \frac{|a_{11}| [V(b_2) - (a_{21}/a_{11}) \text{Cov}(b_1, b_2)]}{|a_{21}| [V(b_1) - (a_{11}/a_{21}) \text{Cov}(b_1, b_2)]}.$$

Now $p_0 \geq 1$ if and only if $\theta \geq 1$, in which case

$$v(p_0) = \frac{V(b_1)V(b_2) - \text{Cov}(b_1, b_2)}{a_{11}^2 V(b_2) + a_{21}^2 V(b_1) - 2a_{11}a_{21} \text{Cov}(b_1, b_2)}.$$

If $\theta \geq 1$, we choose \hat{p} such that $v(\hat{p}) = \min\{v(1), v(p_0), v(\infty)\}$; otherwise, we choose \hat{p} such that $v(\hat{p}) = \min\{v(1), v(\infty)\}$. If, in addition, b_1 and b_2 are uncorrelated, then $\theta = (|a_{11}| V(b_2))/(|a_{21}| V(b_1))$ and we find explicitly that

$$\begin{aligned} \hat{p} &= \infty & 0 < \theta \leq 1 \\ &= 1 + \frac{\ln(|a_{11}|/|a_{21}|)}{\ln|\theta|} & \theta > 1. \end{aligned}$$

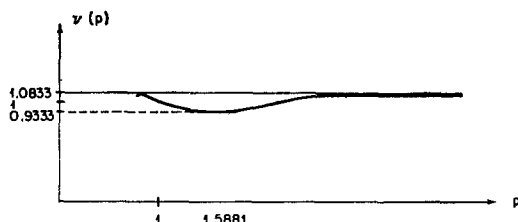
We now present two examples, the solutions of which were obtained utilizing a computer. $v(p)$ was calculated according to (3.3), thus deciding \hat{p} . The corresponding graph was also provided. Note that any type of distribution with $V(b_i) < \infty$ for $i = 1, 2, \dots, n+1$ and yielding the same first two moments would produce the same choice of \hat{p} .

EXAMPLE 1. $A = \begin{pmatrix} 2 \\ 1 \end{pmatrix}$ and the joint distribution of b is

$b_1 \backslash b_2$	1	5
2	$\frac{3}{8}$	$\frac{1}{8}$
6	$\frac{1}{4}$	$\frac{1}{4}$

Hence $E(b_1) = 4$, $E(b_2) = 2.5$, $V(b_1) = 4$, $V(b_2) = 3.75$, $\text{Cov}(b_1, b_2) = 1$. Here $\theta > 1$ and therefore (3.4) yields $p_0 = 1.5881$. Since $v(1) = 1$, $v(p_0) = 0.9333$ and $v(\infty) = 1.0833$, we choose $\hat{p} = 1.5881$. Hence our preferred approximate solution is

$$E(\xi_{1.5881}) = 2.0666,$$



EXAMPLE 2.

$$A = \begin{pmatrix} 0.5 & 1 & 0.5 \\ -3 & -9 & -6 \\ 4 & 12 & 6 \\ 2.5 & 6.5 & 3.5 \end{pmatrix}$$

and the joint distribution of b is $b_4 = 0$,

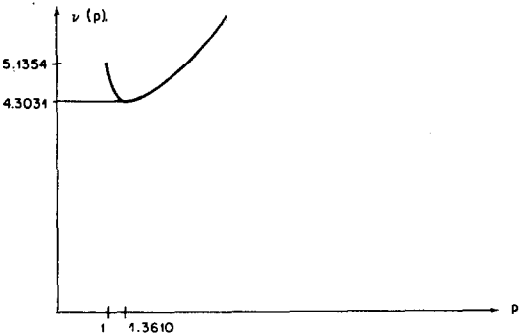
$b_3 = 4$			$b_3 = 6$		
$b_1 \backslash b_2$	1	5	$b_1 \backslash b_2$	1	5
2	1/8	1/4	2	1/16	1/16
6	1/8	0	6	1/8	1/4

Thus

$E(b_1) = 4, \quad E(b_2) = 3.25, \quad E(b_3) = 5, \quad E(b_4) = 0,$
 $V(b_1) = 4, \quad V(b_2) = 3.9375, \quad V(b_3) = 1, \quad V(b_4) = 0,$
 $\text{Cov}(b_1, b_2) = -0.5, \quad \text{Cov}(b_1, b_3) = 1, \quad \text{Cov}(b_2, b_3) = 0.25,$
 $\text{Cov}(b_1, b_4) = \text{Cov}(b_2, b_4) = \text{Cov}(b_3, b_4) = 0.$

Here $\hat{p} = 1.3610$,

$$E(\xi_{1.3610}) = \begin{pmatrix} -2.8273 \\ 3.6896 \\ -4.6631 \end{pmatrix},$$



APPENDIX: SOME RESULTS ON MATRICES

As a consequence of the Binet–Cauchy theorem [3], we have

LEMMA A.1. *If A is an $(n + 1) \times n$ matrix, then $\det(A^T A) = \Delta_2$. We note that for $A \in \mathcal{H}$, $\Delta_2 \neq 0$.*

THEOREM A.1. *If $A \in \mathcal{H}$, then*

$$\begin{aligned} I - A(A^T A)^{-1} A^T \\ &= \frac{1}{\Delta_2} \begin{pmatrix} D_1^2 & -D_1 D_2 & \cdots & (-1)^n D_1 D_{n+1} \\ -D_2 D_1 & D_2^2 & \cdots & (-1)^{n+1} D_2 D_{n+1} \\ \vdots & \vdots & \ddots & \vdots \\ (-1)^n D_{n+1} D_1 & (-1)^{n+1} D_{n+1} D_2 & \cdots & (-1)^{2n} D_{n+1}^2 \end{pmatrix} \\ &= \frac{\delta \delta^T}{\Delta_2}. \end{aligned}$$

Proof. Since $AP = G$, it follows that

$$(A^T A)^{-1} A^T = P(G^T G)^{-1} G^T \quad (\text{A.1})$$

and

$$A(A^T A)^{-1} A^T = G(G^T G)^{-1} G^T. \quad (\text{A.2})$$

From the definition of G^k , we see that

$$\begin{aligned} \hat{D}_k &= (-1)^{n-k} h_k \prod_{i \neq k} g_i, \quad k = 1, 2, \dots, n, \\ &= \prod_{i=1}^n g_i, \quad k = n + 1. \end{aligned}$$

Since $\hat{D}_k = D_k \det(P)$ and $|\det(P)| = 1$, we have $D_i D_j = \hat{D}_i \hat{D}_j$ for $i, j = 1, 2, \dots, n + 1$. From Lemma A.1,

$$\det(A^T A) = \sum_{k=1}^{n+1} D_k^2 = \sum_{k=1}^{n+1} \hat{D}_k^2 = \sum_{k=1}^n h_k^2 \prod_{i \neq k} g_i^2 + \prod_{i=1}^n g_i^2.$$

Since

$$G^T G = \begin{pmatrix} g_1^2 + h_1^2 & h_1 h_2 & \cdots & h_1 h_n \\ h_2 h_1 & g_2^2 + h_2^2 & \cdots & h_2 h_n \\ \vdots & \vdots & \ddots & \vdots \\ h_n h_1 & h_n h_2 & \cdots & g_n^2 + h_n^2 \end{pmatrix},$$

we have

$$(G^T G)^{-1} = \frac{1}{A_2} \begin{bmatrix} \left(\sum_{k \neq 1} \frac{h_k^2}{g_n^2} + 1 \right) \prod_{i \neq 1} g_i^2 & -h_1 h_2 \prod_{i \neq 1, 2} g_i^2 & \cdots & -h_1 h_n \prod_{i \neq 1, n} g_i^2 \\ -h_1 h_2 \prod_{i \neq 2, 1} g_i^2 & \left(\sum_{k \neq 2} \frac{h_k^2}{g_n^2} + 1 \right) \prod_{i \neq 2} g_i^2 & \cdots & -h_2 h_n \prod_{i \neq 2, n} g_i^2 \\ \vdots & \vdots & \ddots & \vdots \\ -h_n h_1 \prod_{i \neq n, 1} g_i^2 & -h_n h_2 \prod_{i \neq n, 2} g_i^2 & \cdots & \left(\sum_{k \neq n} \frac{h_k^2}{g_n^2} + 1 \right) \prod_{i \neq n} g_i^2 \end{bmatrix}.$$

Hence

$$G(G^T G)^{-1} G^T = \frac{1}{A_2} \begin{bmatrix} \sum_{k \neq 1} D_k^2 & D_1 D_2 & -D_1 D_3 & \cdots & (-1)^{n+1} D_1 D_{n+1} \\ D_2 D_1 & \sum_{k \neq 2} D_k^2 & D_2 D_3 & \cdots & (-1)^{n+2} D_2 D_{n+1} \\ \vdots & \vdots & \vdots & \ddots & \vdots \\ (-1)^{n+1} D_{n+1} D_1 & (-1)^{n+2} D_{n+1} D_2 & (-1)^{n+3} D_{n+1} D_3 & \cdots & \sum_{k \neq n+1} D_k^2 \end{bmatrix}.$$

The result now follows in view of Lemma A.1 and (A.2).

LEMMA A.2. If $A \in \mathcal{H}$ and M_1 and M_2 are both $n \times r$ matrices then $AM_1 = AM_2$ implies that $M_1 = M_2$.

THEOREM A.2. If $A \in \mathcal{H}$, then $(A^T A)^{-1} A^T = M = (m_{ij})$, where m_{ij} is defined in (2.6).

Proof. By Lemma A.2, this theorem is proved if $AM = A(A^T A)^{-1} A^T$. The entry of AM at the l th row and the j th column is

$$\sum_{i=1}^n a_{li} m_{ij} = \frac{1}{A_2} \begin{vmatrix} a_{11} & a_{12} & \cdots & a_{1n} & (-1)^n D_1 \\ \vdots & \vdots & & \vdots & \vdots \\ a_{j-1,1} & a_{j-1,2} & \cdots & a_{j-1,n} & (-1)^{n-j+2} D_{j-1} \\ a_{l,1} & a_{l,2} & \cdots & a_{l,n} & 0 \\ a_{j+1,1} & a_{j+1,2} & \cdots & a_{j+1,n} & (-1)^{n-1} D_{j+1} \\ \vdots & \vdots & \cdots & \vdots & \vdots \\ a_{n+1,1} & a_{n+1,2} & \cdots & a_{n+1,n} & (-1)^0 D_{n+1} \end{vmatrix},$$

where $l = 1, \dots, n+1, j = 1, \dots, n+1$.

If $l = j$, the determinant expanded about the last column yields

$$\sum_{i=1}^n a_{ji} m_{ij} = (1/\Delta_2) \sum_{i \neq j} D_i^2.$$

If $l \neq j$, the same expansion yields zero for each entry except the l th one, for which we have

$$\sum_{i=1}^n a_{li} m_{ij} = ((-1)^{l+j-1}/\Delta_2) D_l D_j.$$

Therefore, we see that matrix AM is precisely $G(G^T G)^{-1} G^T = A(A^T A)^{-1} A^T$ as constructed in Theorem A.1.

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